BLOCK-DIAGONAL AND CONSTRAINT PRECONDITIONERS FOR NONSYMMETRIC INDEFINITE LINEAR SYSTEMS. PART I: THEORY

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Abstract. We study block diagonal preconditioners and an efficient variant of constraint preconditioners for general two-by-two block linear systems with zero (2,2) block. We derive block diagonal preconditioners from a splitting of the (1,1)-block of the matrix. From the resulting preconditioned system we derive a smaller, so-called 'related' system that yields the solution of the original problem. Solving the related system corresponds to an efficient implementation of constraint preconditioning. We analyze the properties of both classes of preconditioned matrices, in particular their spectrum. Using analytical results we show that the related system matrix has the more favorable spectrum, which in many applications translates into faster convergence for Krylov subspace methods. We show that fast convergence depends mainly on the quality of the splitting, a topic for which a substantial body of theory exists. Our analysis also provides a number of new relations between block-diagonal preconditioners and constraint preconditioners. For constrained problems, solving the the related system produces iterates that satisfy the constraints exactly, just as for systems with a constraint preconditioner. Finally, for the Lagrange multiplier formulation of a constrained optimization problem we show how scaling nonlinear constraints can dramatically improve the convergence for linear systems in a Newton iteration. Our theoretical results are confirmed by numerical experiments on a constrained optimization problem.

Our approach is very general, as we make almost no assumptions on the given block two-by-two system and the splitting that defines the preconditioners. In particular, the system matrix might be nonsymmetric, and the (1,1) block might be indefinite, or even singular. This is the first paper in a two-part sequence. In the second paper we will study the use of our preconditioners in a variety of applications.

Key words. Saddle point systems, indefinite systems, eigenvalue bounds, Krylov subspace methods, preconditioning, constrained optimization, mesh-flattening

AMS subject classifications. 65F10, 65F15, 65D18

1. Introduction.

We study preconditioners for general nonsingular linear systems of the type

$$\mathcal{A}u = \begin{bmatrix} A & B^T \\ C & 0 \end{bmatrix} \begin{bmatrix} \tilde{x} \\ \tilde{y} \end{bmatrix} = \begin{bmatrix} \tilde{f} \\ \tilde{g} \end{bmatrix}, \qquad A \in \mathbb{R}^{n\times n}, \quad B, C \in \mathbb{R}^{m\times n}, \text{ with } n \geq m. \tag{1.1}$$

Such systems arise in a large number of applications, for example, the (linearized) Navier-Stokes equations and other physical problems with conservation laws as well as constrained optimization problems.

As such systems are typically large and sparse, solution by iterative methods has been studied extensively. Much attention has focused on the Navier-Stokes problem; see, e.g., [6, 14, 28, 29, 32]. The techniques for solving systems like (1.1) are

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so numerous that it is almost impossible to give an overview. In addition to the methods developed specifically for Navier-Stokes problems, existing techniques also include Uzawa-type algorithms [7], splitting schemes [4, 10], constraint preconditioning [10, 19, 21, 22, 15, 23], and (preconditioned) Krylov subspace methods based on (approximations to) the Schur complement [1, 16, 20].

We start with block-diagonal preconditioners for the general system (1.1); see Section 2 for our assumptions. Results for the general system have been obtained before, for example, Murphy et al. [20] propose the block-diagonal preconditioner

$$\begin{bmatrix} A^{-1} & 0 \\ 0 & (CA^{-1}B^T)^{-1} \end{bmatrix}.$$

If defined, this preconditioner leads to (left or right) preconditioned matrices that are diagonalizable and have at most three eigenvalues. Hence, a Krylov subspace method with the optimality or the Galerkin property, e.g., GMRES [24] or BiCG [8], will converge in at most three steps [20, Remark 3]. However, this preconditioner is more expensive than direct solution by block elimination. So, one typically uses approximations to A^{-1} and $(CA^{-1}B^T)^{-1}$.

We derive such approximations from a splitting of the (1,1) block, $A = D - E$, where D can be efficiently inverted. Then, from a splitting of the preconditioned matrix we derive a fixed point iteration and its so-called 'related' system [13] that have (significantly) fewer unknowns. We show that the related system is not only smaller than the preconditioned system, but it typically also leads to faster convergence for the GMRES algorithm. Furthermore, as solving the related system corresponds to an efficient implementation of constraint preconditioning, each GMRES iterate for the related system satisfies the constraints in (1.1) exactly. For many applications this is important and may lead to significant savings. For a special starting guess a similar result is given in [21], but the potential savings of solving a much smaller system are not elaborated. To analyze the preconditioned iterations derived from (1.1) we introduce a new decomposition to analyze oblique projections using the (cosines of) principal angles between the range and null space of the projector (see (4.10)–(4.11)). We believe this decomposition, which is derived from the SVD that gives the principal angles between two spaces [9, Section 12.4], is useful more generally where the detailed relations between two spaces play a role. Finally, for the Lagrange multiplier formulation of a constrained optimization problem we show how scaling nonlinear constraints can dramatically improve the convergence for linear systems in a Newton iteration.

This paper is the first of a two-part sequence. Here, we focus on the derivation and analysis of our preconditioned iterations, and succinctly demonstrate our results numerically. In part II we demonstrate our results on a variety of applications and discuss efficient implementations of our preconditioners [17].

The paper is organized as follows. In Section 2 we introduce the block-diagonal preconditioners and the resulting preconditioned systems. In Section 3 we study the properties of the preconditioned matrices, in particular their eigendecompositions. In Section 4 we derive the fixed point iteration and its related system. We also analyze the spectral radius of the fixed point iteration matrix and the spectrum of the related system matrix. In Section 5 we discuss why typically it is more efficient to apply GMRES to the related system than to the preconditioned system. In Section 6 we introduce the application for our numerical experiments, and we discuss the scaling of nonlinear constraints in this constrained optimization problem to improve the con-

vergence in each Newton step. In Section 7 we show a few numerical results, and in Section 8 we give our conclusions.

2. Block-diagonal preconditioners. Suppose that a system of the form (1.1) is given, and that we split the (1,1) block of $\mathcal{A}$ into

$$A = D - E, \tag{2.1}$$

where D is invertible. Our *only assumptions* on $\mathcal{A}$ are that this matrix is invertible and the invertible matrix D can be chosen so that $CD^{-1}B^T$ is invertible as well.

Since D and $CD^{-1}B^T$ are both invertible, we can define the preconditioner

$$\mathcal{P}(D) = \begin{bmatrix} D^{-1} & 0 \\ 0 & (CD^{-1}B^T)^{-1} \end{bmatrix}. \tag{2.2}$$

Multiplying $\mathcal{A}$ from the left or right by $\mathcal{P}(D)$ results in the matrices

$$\begin{bmatrix} I_n - D^{-1}E & D^{-1}B^T \\ (CD^{-1}B^T)^{-1}C & 0 \end{bmatrix} \text{ or } \begin{bmatrix} I_n - ED^{-1} & B^T(CD^{-1}B^T)^{-1} \\ CD^{-1} & 0 \end{bmatrix}, \tag{2.3}$$

respectively. Both of these matrices are of the form

$$\mathcal{B}(S) = \begin{bmatrix} I_n - S & N \\ M & 0 \end{bmatrix}, \text{ where} \tag{2.4}$$

$$MN = I_m, \ (NM)^2 = NM, \ S \in \mathbb{R}^{n\times n}, \ M, N^T \in \mathbb{R}^{m\times n}, \ n \geq m. \tag{2.5}$$

After applying $\mathcal{P}(D)$ to (1.1) we are interested in solving linear systems of the form

$$\mathcal{B}(S)\begin{bmatrix} x \\ y \end{bmatrix} = \begin{bmatrix} f \\ g \end{bmatrix}. \tag{2.6}$$

Note that either the vector $[\tilde{x}^T, \tilde{y}^T]^T$ or $[\tilde{f}^T, \tilde{g}^T]^T$ from the original problem (1.1) is modified to account for the application of either right or left preconditioning.

REMARK 2.1. If $n = m$, the matrix $CD^{-1}B^T$ is invertible if and only if both C and B^T are invertible. In this case we can solve (1.1) directly by computing $x = C^{-1}g$ and $y = B^{-T}(f - Ax)$. This has essentially the same cost as one multiplication of (1.1) by (2.2), and preconditioning has no advantage over solving (1.1) directly. While many of our results hold true for $n = m$, we consider this case of little interest.

Our general approach is to consider which splittings $A = D - E$ result in preconditioned systems (2.6) that are solved efficiently by an iterative method. If we consider only the iteration count, the most effective preconditioner of the form (2.2) was derived by Murphy et al. [20]. In our notation it is $\mathcal{P}(A)$, corresponding to the trivial splitting $D = A$ and $E = 0$. As shown in [20], the left and right preconditioned matrices, both of the form $\mathcal{B}(0)$, are diagonalizable with at most three distinct eigenvalues in the nonsingular case. Hence, any Krylov subspace method with a Galerkin or optimality property, e.g. BiCG [8] or GMRES [24] (see [12] for an overview), will converge in at most three steps. While this is an attractive feature, the preconditioner $\mathcal{P}(A)$ requires multiplications by A^{-1}. However, in many applications we have $n \gg m$, and the computational effort to solve for A is not significantly less than the

effort to solve for $\mathcal{A}$. In addition, a difficult subproblem arises in inverting the Schur complement $CA^{-1}B^T$. Since A^{-1} is usually not available explicitly, $CA^{-1}B^T$ cannot be formed without solving for A a significant number of times. Murphy et al. were, of course, aware that using $\mathcal{P}(A)$ typically is prohibitive, and they remarked that approximations to A^{-1} and $(CA^{-1}B^T)^{-1}$ should lead to clustered eigenvalues as well, where the clustering "depends on the quality of the approximations" [20, Remark 5]. To some extent this remark was the motivation for part of our work.

Summarizing, the general strategy must be to choose a splitting that leads to efficiently invertible matrices D and $CD^{-1}B^T$ and preserves the properties of the algebraically optimal preconditioner $\mathcal{P}(A)$ as much as possible. To derive guidelines for such choices we analyze how properties of the preconditioned matrices depend on the splitting.

3. Properties of the Matrices $\mathcal{B}(S)$. Our first goal is to identify the conditions under which $\mathcal{B}(S)$ is singular.

THEOREM 3.1. *A matrix $\mathcal{B}(S)$ of the form (2.4)–(2.5) is singular if, and only if, 1 (one) is an eigenvalue of the matrix $(I_n - NM)S$. In particular, each matrix $\mathcal{B}(0)$ of the form (2.4)–(2.5) is nonsingular.*

Proof. The matrix $\mathcal{B}(S)$ is singular if, and only if, there exists a nonzero vector $[x^T, y^T]^T$ for which $\mathcal{B}(S)\,[x^T, y^T]^T = 0$. This is equivalent to the two equations

$$(i)\;\; (I_n - S)x + Ny \;=\; 0, \quad (ii)\;\; Mx \;=\; 0.$$

Equation (i) is equivalent to $Sx - Ny = x$. Inserting this into equation (ii), and using that $MN = I_m$, yields $y = MSx$. Hence $x = 0$ implies $y = 0$. Inserting $y = MSx$ into $Sx - Ny = x$ shows that x has to satisfy $(I_n - NM)Sx = x$. There exists a nonzero x satisfying this requirement, i.e., $\mathcal{B}(S)$ is singular, if, and only if, $(I_n - NM)S$ has an eigenvalue 1. ∎

REMARK 3.2. Under the assumption that the preconditioner $\mathcal{P}(A)$ exists, i.e., A and $CA^{-1}B^T$ are both invertible, Theorem 3.1 shows that the matrices $\mathcal{P}(A)\,\mathcal{A}$ and $\mathcal{A}\,\mathcal{P}(A)$ are always nonsingular. Hence the zero eigenvalue included in the discussion of Murphy et al. (cf. [20, Remark 1]) never occurs.

The inverse of $\mathcal{B}(0)$ can be easily computed and is given by

$$\mathcal{B}(0)^{-1} \;=\; \begin{bmatrix} I_n - NM & N \\ M & -I_m \end{bmatrix} \;=\; \mathcal{B}(0) - \begin{bmatrix} NM & 0 \\ 0 & I_m \end{bmatrix}. \tag{3.1}$$

In case of the trivial splitting ($S = 0$), one can therefore simply solve (2.6) via (3.1), which yields

$$(i)\;\; x = (I_n - NM)f + Ng, \quad (ii)\;\; y = Mf - g.$$

The solution is computed using two matrix-vector products with N, one matrix-vector product with M, and three vector additions. Therefore, the cost of this solution method is comparable to that of just one step of a Krylov subspace method applied to (2.6).

Next, we study the properties of eigenvalues and eigenvectors of $\mathcal{B}(S)$. In the case $S = 0$, we can directly relate the eigenvalues and eigenvectors of $\mathcal{B}(0)$ to the projection matrix NM.

THEOREM 3.3. *Let $\mathcal{B}(0)$ be of the form (2.4)–(2.5). Then $\mathcal{B}(0)$ is diagonalizable, and it has*

- *$n-m$ eigenpairs of the form $(1,\ [u_j^T,\ 0]^T)$, where $u_1,\ldots,u_{n-m}$ form a basis of $Null(NM)$, the nullspace of NM.*
- *$2m$ eigenpairs of the form $(\lambda^{\pm},\ [u_j^T,\ (\lambda^{\pm})^{-1}(Mu_j)^T]^T)$, where $\lambda^{\pm} \equiv (1 \pm \sqrt{5})/2$, and $u_{n-m+1},\ldots,u_n$ form a basis of $Range(NM)$, the range of NM.*

In particular, if we denote

$$U_1 \ \equiv\ [u_1,\ldots,u_{n-m}] \ \in \mathbb{R}^{n\times(n-m)},\ \ U_2 \ \equiv\ [u_{n-m+1},\ldots,u_n] \ \in \mathbb{R}^{n\times m}\ , \qquad (3.2)$$

then the eigenvector matrix $\mathcal{Y}(0)$ of $\mathcal{B}(0)$ is given by

$$\mathcal{Y}(0) \ =\ \begin{bmatrix} U_1 & U_2 & U_2 \\ 0 & (\lambda^+)^{-1}MU_2 & (\lambda^-)^{-1}MU_2 \end{bmatrix}, \qquad (3.3)$$

where both $[U_1, U_2] \in \mathbb{R}^{n\times n}$ and $(\lambda^-)^{-1}MU_2 \ \in \mathbb{R}^{m\times m}$ are nonsingular.

Proof. To compute the eigendecomposition of $\mathcal{B}(0)$, we consider the equation $\mathcal{B}(0)\,[u^T,\,v^T]^T \ =\ \lambda\,[u^T,\,v^T]^T$, which is equivalent to the two equations

$$(i)\ u + Nv \ =\ \lambda u, \quad (ii)\ Mu \ =\ \lambda v.$$

Since $\mathcal{B}(0)$ is nonsingular we can assume that $\lambda \neq 0$, so that equation (ii) is equivalent to $v = \lambda^{-1}Mu$. Inserting this into (i), and multiplying the resulting equation by λ yields

$$NMu = (\lambda^2 - \lambda)u,$$

i.e., the u-component of each eigenvector $[u^T,\,v^T]^T$ of $\mathcal{B}(0)$ is an eigenvector of the projection NM. Hence $\lambda^2 - \lambda$ is either equal to one (i.e. $\lambda = (1 \pm \sqrt{5})/2$), or equal to zero (i.e. $\lambda = 1$).

Next note that since $MN = I_m$, we have

$$m \ =\ Rank(MN) \ \leq\ \min\left(Rank(M),\ Rank(N)\right) \ \leq\ m.$$

Hence $Rank(N) = m$, so that $Null(NM)$ is equal to $Null(M)$ and has dimension $n-m$. If $u_1,\ldots,u_{n-m}$ form a basis of $Null(NM)$, then the $n-m$ pairs $(1,\ [u_j^T,\ 0]^T)$, $j = 1,\ldots,n-m$, satisfy equations (i) and (ii). Furthermore, let the m vectors $u_{n-m+1},\ldots,u_n$ form a basis of $Range(NM)$. Using these vectors in (i) and (ii) shows that the remaining $2m$ eigenpairs are $(\lambda^{\pm},\ [u_j^T,\ (\lambda^{\pm})^{-1}(Mu_j)^T]^T)$, $j = n-m+1,\ldots,n$, with $\lambda^{\pm} \equiv (1 \pm \sqrt{5})/2$.

Finally, $[U_1, U_2]$ is nonsingular since this matrix is the eigenvector matrix of the projection NM. Furthermore, if MU_2 were singular, then a nonzero vector w would exist such that $MU_2w = 0$. However, multiplication with N yields $NMU_2w = U_2w = 0$, which is a contradiction since the columns of U_2 are linearly independent. ∎

REMARK 3.4. The statement of Theorem 3.3 contains the complete eigendecompositions of the preconditioned matrices $\mathcal{P}(A)\,\mathcal{A}$ and $\mathcal{A}\,\mathcal{P}(A)$ that are the subject of [20]. In that paper Murphy et al. show that the two matrices are diagonalizable and derive the location of their eigenvalues.

Note that in the case $n = m$, we have $Null(NM) = \{0\}$, and Theorem 3.3 shows that in this case the only eigenvalues of $\mathcal{B}(0)$ are $(1 \pm \sqrt{5})/2$. As discussed in Remark 2.1, this case is of little interest for our purposes. *In the following we will therefore assume that $n > m$.*

Next, we derive bounds on the eigenvalues of each matrix $\mathcal{B}(S)$ in terms of the corresponding matrix $\mathcal{B}(0)$.

THEOREM 3.5. *Consider matrices $\mathcal{B}(S)$ of the form (2.4)–(2.5) with fixed N and M. Let $\mathcal{B}(0) = \mathcal{Y}(0)\,\mathcal{D}\,\mathcal{Y}(0)^{-1}$ denote the eigendecomposition of $\mathcal{B}(0)$, with eigenvector matrix $\mathcal{Y}(0)$ given as in (3.3), and let $[U_1, U_2]$ denote the corresponding eigenvector matrix of the projection NM. Then for each matrix S, and each eigenvalue λ_S of $\mathcal{B}(S)$, there is an eigenvalue λ of $\mathcal{B}(0)$ such that*

$$|\lambda_S - \lambda| \;\le\; \left\| \mathcal{Y}(0)^{-1} \begin{bmatrix} S & 0 \\ 0 & 0 \end{bmatrix} \mathcal{Y}(0) \right\| \tag{3.4}$$

$$\le\; c_S \,\| \,[U_1, U_2]^{-1}\, S\, [U_1, U_2] \,\|\,, \tag{3.5}$$

where $c_S \equiv \left(\, 2 + \frac{2}{5}(\lambda^-)^2 \,\right)^{\frac{1}{2}} \approx 1.4672$.

Proof. A given matrix $\mathcal{B}(S)$ can be additively split into

$$\mathcal{B}(S) \;=\; \mathcal{B}(0) - \begin{bmatrix} S & 0 \\ 0 & 0 \end{bmatrix}. \tag{3.6}$$

Since $\mathcal{B}(0)$ is diagonalizable, and has eigenvector matrix $\mathcal{Y}(0)$, inequality (3.4) follows from a well-known result in matrix perturbation theory [30, Theorem IV.1.12].

To prove (3.5) consider the following two-by-two block decomposition

$$\mathcal{Y}(0) \;=\; \begin{bmatrix} Y_{11} & Y_{12} \\ Y_{21} & Y_{22} \end{bmatrix},$$

where $Y_{11} \equiv [U_1, U_2] \in \mathbb{R}^{n\times n}$ and $Y_{22} \equiv (\lambda^-)^{-1} M U_2 \;\in \mathbb{R}^{m\times m}$ are both invertible, cf. (3.3). Then the inverse of $\mathcal{Y}(0)$ satisfies

$$\mathcal{Y}(0)^{-1} = \begin{bmatrix} (Y_{11} - Y_{12}Y_{22}^{-1}Y_{21})^{-1} & -Y_{11}^{-1}Y_{12}\,(Y_{22} - Y_{21}Y_{11}^{-1}Y_{12})^{-1} \\ -Y_{22}^{-1}Y_{21}\,(Y_{11} - Y_{12}Y_{22}^{-1}Y_{21})^{-1} & (Y_{22} - Y_{21}Y_{11}^{-1}Y_{12})^{-1} \end{bmatrix}.$$

An elementary computation now shows that

$$\begin{aligned}
(Y_{11} - Y_{12}Y_{22}^{-1}Y_{21})^{-1} &= \big(\,[U_1, U_2] - U_2(\lambda^-)(MU_2)^{-1}\,[0, (\lambda^+)^{-1}MU_2]\,\big)^{-1} \\
&= \big(\,[U_1, U_2] - [0, (\lambda^-/\lambda^+)U_2]\,\big)^{-1} \\
&= \begin{bmatrix} I_{n-m} & 0 \\ 0 & (\lambda^+/\sqrt{5})I_m \end{bmatrix} [U_1, U_2]^{-1} \\
&\equiv \hat{I}_n\,[U_1, U_2]^{-1}, \text{ and}
\end{aligned}$$

$$-Y_{22}^{-1}Y_{21}\,(Y_{11} - Y_{12}Y_{22}^{-1}Y_{21})^{-1} = -[0,\, (\lambda^-/\sqrt{5})I_m]\,[U_1, U_2]^{-1} \;\equiv\; \hat{I}_m\,[U_1, U_2]^{-1}.$$

Using these relations the square of the right hand side of (3.4) is equal to

$$\begin{aligned}
&\left\| \begin{bmatrix} \hat{I}_n[U_1, U_2]^{-1} S [U_1, U_2] & \hat{I}_n[U_1, U_2]^{-1} S U_2 \\ \hat{I}_m[U_1, U_2]^{-1} S [U_1, U_2] & \hat{I}_m[U_1, U_2]^{-1} S U_2 \end{bmatrix} \right\|^2 \\
&= \max_{\|[a,b]\|=1} \left\| \begin{bmatrix} \hat{I}_n[U_1, U_2]^{-1} S\,([U_1, U_2]a + U_2 b) \\ \hat{I}_m[U_1, U_2]^{-1} S\,([U_1, U_2]a + U_2 b) \end{bmatrix} \right\|^2 \\
&= \max_{\|[a_1,a_2,b]\|=1} \left\| \begin{bmatrix} \hat{I}_n[U_1, U_2]^{-1} S\,(U_1 a_1 + U_2(a_2 + b)) \\ \hat{I}_m[U_1, U_2]^{-1} S\,(U_1 a_1 + U_2(a_2 + b)) \end{bmatrix} \right\|^2 \\
&\le \max_{\|c\|\le\sqrt{2}} \left\| \begin{bmatrix} \hat{I}_n[U_1, U_2]^{-1} S\,[U_1, U_2] c \\ \hat{I}_m[U_1, U_2]^{-1} S\,[U_1, U_2] c \end{bmatrix} \right\|^2
\end{aligned}$$

$$\leq 2\left(\|\hat{I}_n[U_1,U_2]^{-1}S[U_1,U_2]\|^2 + \|\hat{I}_m[U_1,U_2]^{-1}S[U_1,U_2]\|^2\right)$$
$$\leq 2\left(1+(\lambda^-/\sqrt{5})^2\right)\ \|[U_1,U_2]^{-1}S[U_1,U_2]\|^2.$$

Taking square roots completes the proof. ∎

Each choice of the splitting $A = D - E$ leads to *fixed* matrices $\mathcal{B}(0)$ and $\mathcal{B}(S)$ for which Theorem 3.5 will hold. Hence, the theorem, which allows S to vary, is more general than our application requires. Furthermore, in bounding the right hand side of (3.4) from above by (3.5) we have used three inequalities, which should generally be tight. So, we can expect the right hand side of (3.4) to be close to (3.5).

To analyze implications of Theorem 3.5, recall that the columns of U_1 and U_2 form the bases of $Null(NM)$ and $Range(NM)$, respectively. We choose both bases to be orthonormal, i.e., $U_1^T U_1 = I_{n-m}$ and $U_2^T U_2 = I_m$. A key ingredient of our analysis is the singular value decomposition (SVD) of the matrix $U_1^T U_2$,

$$U_1^T\, U_2 \;=\; \Phi\Omega\Psi^T \;=\; [\varphi_1,\ldots,\varphi_{n-m}]\,\mathrm{diag}(\omega_1,\ldots,\omega_k)\,[\psi_1,\ldots,\psi_m]^T, \tag{3.7}$$

where $\Phi \in \mathbb{R}^{(n-m)\times(n-m)}$ and $\Psi \in \mathbb{R}^{m\times m}$ are both orthogonal matrices, and $\Omega \in \mathbb{R}^{(n-m)\times m}$ with $\omega_1 \geq \omega_2 \geq \cdots \geq \omega_k$, and $k = \min(n-m, m)$. It is well known that the singular values satisfy $\omega_j = \cos(\theta_j)$, where the θ_j are the principal angles between $Null(NM)$ and $Range(NM)$; see [9, Section 12.4]. Since NM is a projection, we have $Null(NM) \cap Range(NM) = \{0\}$, and thus $\omega_j \in [0,1)$, $j = 1,\ldots,k$.

LEMMA 3.6. *Let $U_1^T U_1 = I_{n-m}$, $U_2^T U_2 = I_m$, and let the SVD (3.7) be defined. Then $[U_1,U_2]$ has $2k$ singular values $(1\pm\omega_j)^{1/2}$, $j = 1,2,\ldots,k$, and an $(n-2k)$-fold singular value 1 (one), where $k = \min(n-m,m)$. In particular, the condition number of $[U_1,U_2]$ is given by*

$$\kappa([U_1,U_2]) \;=\; \left(\frac{1+\omega_1}{1-\omega_1}\right)^{1/2}. \tag{3.8}$$

Proof. Since

$$[U_1,U_2]^T\,[U_1,U_2] \;=\; \begin{bmatrix} \Phi & 0 \\ 0 & \Psi \end{bmatrix}\begin{bmatrix} I_{n-m} & \Omega \\ \Omega^T & I_m \end{bmatrix}\begin{bmatrix} \Phi^T & 0 \\ 0 & \Psi^T \end{bmatrix},$$

the singular values of $[U_1,U_2]$ are the square roots of the eigenvalues of

$$\begin{bmatrix} I_{n-m} & \Omega \\ \Omega^T & I_m \end{bmatrix}. \tag{3.9}$$

$(\lambda, [z^T, v^T]^T)$ is an eigenpair of this matrix if and only if

$$(i)\ z+\Omega v \;=\; \lambda z, \quad (ii)\ \Omega^T z + v \;=\; \lambda v.$$

Let us assume that $n-m \geq m$. From (ii) we have $\Omega^T z = (\lambda-1)v \Rightarrow \Omega\Omega^T z = (\lambda-1)\Omega v$, and (i) gives $\Omega v = (\lambda-1)z$. Substituting for Ωv into the former equation gives

$$\Omega\Omega^T z = (\lambda-1)^2 z\,. \tag{3.10}$$

First, assume that $\omega_k \neq 0$ for $k = 1, \ldots, m$. Then $\Omega\Omega^T$ is a diagonal matrix with the m leading nonzero coefficients $\omega_1^2, \ldots, \omega_m^2$, followed by $n - 2m$ zeros, and we have the following solutions of (3.10), which give the eigenpairs of (3.9).

For $k = 1, \ldots, m$ we can take $z = e_k$ and solve $\omega_k^2 = (\lambda - 1)^2$. This gives $\omega_k = \pm(\lambda - 1) \Rightarrow \lambda = 1 \pm \omega_k$. Clearly, when ω_k is not unique we have additional degrees of freedom. From (i), we see that for the eigenvalues $\lambda = 1 \pm \omega_k$ we find the vectors $v = \pm e_k$. Hence, we have the two eigenpairs $\{1 + \omega_k, [e_k^T, e_k^T]^T\}$ and $\{1 - \omega_k, [e_k^T, -e_k^T]^T\}$ for $k = 1, \ldots, m$.

For $k = m + 1, \ldots, n - m$ we can take $z = e_k$ and solve $(\lambda - 1)^2 = 0 \Rightarrow \lambda = 1$. Of course, any basis for $Null(\Omega\Omega^T)$ leads to a valid choice for the vectors z. For all choices, (i) gives $\Omega v = 0$ which implies $v = 0$ since the columns of Ω are independent (assuming $\omega_k \neq 0$ for $k = 1, \ldots, m$). Hence, we have the eigenpairs $\{1, [e_k^T, 0]^T\}$ for $k = m + 1, \ldots, n - m$.

If $\omega_k = 0$ for one or more $k \leq m$, this leads to additional degrees of freedom in choosing the eigenvectors. However, the choices for eigenvalues and eigenvectors above remain correct. Notice that the two eigenvalues $\lambda = 1 \pm \omega_k$ give rise to two additional eigenvalues $\lambda = 1$.

Finally, we find the singular values by taking the (positive) square roots of the eigenvalues λ. From the definition of Ω we see that the largest and smallest singular value are given by respectively $1 + \omega_1$ and $1 - \omega_1$, which proves (3.8). The case $n - m < m$ can be proved analogously. ∎

Lemma 3.6 provides a general result about the eigenvector matrix of a projection. With a little additional algebra its proof also gives the right singular vectors (through the eigenvectors of (3.9)) and the left singular vectors (through multiplication with $[U_1, U_2]$). Using Lemma 3.6 we can simplify the bound on the eigenvalues of $\mathcal{B}(S)$.

COROLLARY 3.7. *In the notation of Theorem 3.5, for each eigenvalue λ_S of $\mathcal{B}(S)$, there is an eigenvalue λ of $\mathcal{B}(0)$, such that*

$$|\lambda_S - \lambda| \;\leq\; c_S \left(\frac{1 + \omega_1}{1 - \omega_1}\right)^{1/2} \|S\|, \tag{3.11}$$

where ω_1 is the largest singular value of $U_1^T U_2$.

In particular, if $\omega_1 = 1 - \varepsilon$, then (3.11) becomes

$$|\lambda_S - \lambda| \;\leq\; c_S \,(2\varepsilon^{-1} - 1)^{1/2} \,\|S\| \approx \sqrt{2}\, c_S \,\varepsilon^{-1/2} \,\|S\| \;\approx\; 2.075\, \varepsilon^{-1/2} \,\|S\| \,. \tag{3.12}$$

Hence, if the angles between $Null(NM)$ and $Range(NM)$ are not too small, the eigenvalue perturbation depends essentially on $\|S\|$. This shows that to derive good preconditioners the main concern is to find a good splitting $A = D - E$. This is a well-developed research area [31, 13].

4. An efficient implementation of the constraint preconditioner. We now derive a smaller system whose solution leads to the solution of the overall system. We extend the so-called constraint preconditioner for the symmetric case [15, 21, 23, 22] to the general case discussed here, and we derive an efficient method of solution.

Rather than apply the constraint preconditioner directly we derive it from the preconditioned system (2.6) to emphasize the relation between the two. Consider the splitting (3.6) derived from (2.6) and the resulting system of linear equations,

$$\mathcal{B}(0) \begin{bmatrix} x \\ y \end{bmatrix} = \begin{bmatrix} S & 0 \\ 0 & 0 \end{bmatrix} \begin{bmatrix} x \\ y \end{bmatrix} + \begin{bmatrix} f \\ g \end{bmatrix}. \tag{4.1}$$

Multiplying both sides from the left by $\mathcal{B}(0)^{-1}$ (see (3.1)) leads to the fixed point iteration

$$\begin{bmatrix} x_{k+1} \\ y_{k+1} \end{bmatrix} = \begin{bmatrix} (I_n - NM)Sx_k \\ MSx_k \end{bmatrix} + \begin{bmatrix} \hat{f} \\ \hat{g} \end{bmatrix}, \tag{4.2}$$

where $[\hat{f}^T, \hat{g}^T]^T \equiv \mathcal{B}(0)^{-1}[f^T, g^T]^T$. Since the right hand side depends only on x_k, the convergence of (4.2) depends only on the iteration

$$x_{k+1} \;=\; (I_n - NM)Sx_k + \hat{f} \;\equiv\; Fx_k + \hat{f}. \tag{4.3}$$

The observation that a splitting of (1.1) based on the constraint preconditioner (4.5) leads to a fixed point iteration that does not depend on y_k is also made in [2]. However, no consequences are mentioned. Note that y_{k+1} in (4.2) is available essentially for free. The iteration (4.3) converges if and only if the spectral radius of $F = (I_n - NM)S$ satisfies $\rho(F) < 1$.

A fixed point x of iteration (4.3) satisfies $x = Fx + \hat{f}$, or, equivalently,

$$Rx \;=\; \hat{f}, \quad R \;\equiv\; I_n - F. \tag{4.4}$$

This is called the 'related' system for the fixed-point iteration [13], and we can solve (4.4) by a Krylov subspace method as an alternative to solving (2.6).

We now turn to the relation between the related system and a constraint preconditioner for the general problem (1.1). In our notation this preconditioner is given by

$$\begin{bmatrix} D & B^T \\ C & 0 \end{bmatrix}^{-1} = \\ \begin{bmatrix} D^{-1} - D^{-1}B^T(CD^{-1}B^T)^{-1}CD^{-1} & -D^{-1}B^T(CD^{-1}B^T)^{-1} \\ (CD^{-1}B^T)^{-1}CD^{-1} & -(CD^{-1}B^T)^{-1} \end{bmatrix}. \tag{4.5}$$

Preconditioning (1.1) from the left by (4.5) yields

$$\begin{bmatrix} D & B^T \\ C & 0 \end{bmatrix}^{-1} \begin{bmatrix} A & B^T \\ C & 0 \end{bmatrix} \begin{bmatrix} \tilde{x} \\ \tilde{y} \end{bmatrix} = \begin{bmatrix} D & B^T \\ C & 0 \end{bmatrix}^{-1} \begin{bmatrix} \tilde{f} \\ \tilde{g} \end{bmatrix} \Leftrightarrow \tag{4.6}$$

$$\begin{bmatrix} I_n - (I_n - NM)S & 0 \\ -MS & I_m \end{bmatrix} \begin{bmatrix} x \\ y \end{bmatrix} = \begin{bmatrix} \hat{f} \\ \hat{g} \end{bmatrix}, \tag{4.7}$$

where, as for the left block-diagonal preconditioned matrix $\mathcal{B}(\mathcal{S})$ given in (2.3), we define $N \equiv D^{-1}B^T$, $M \equiv (CD^{-1}B^T)^{-1}C$, and $S \equiv D^{-1}E$.

It turns out that in case of left (constraint and block-diagonal) preconditioning the (1,1) block of the matrix in (4.7) is precisely the related system matrix R. Hence, solving (4.4) and computing y from (4.2) corresponds to solving the (1,1) block of (4.7) separately and again computing y at the end. Solving the related system rather than the whole system (4.7) has several advantages that are pointed out in Section 5.

Constraint preconditioners, mainly for Hermitian matrices $\mathcal{A}$, have been discussed at many places; most useful for our discussion are [10, 11, 15, 19, 21, 22, 23]. In [11, 15, 19, 21, 22, 23] the matrix is Hermitian, and in [19, 21, 22, 23] the (1,1) block of the matrix either is positive definite or made positive definite by transforming the

problem. In [15] the (1,1) block may be indefinite, but it must be nonsingular. In [10] the (1,1) block may be nonsymmetric, but $B = C$ (in our notation) must hold. In all cases the constraint preconditioner itself is symmetric with a positive definite (1,1) block. In fact, in [19] the (1,1) block is diagonal, and in [21, 22, 23] it is the identity matrix. In [21, 22, 23] the authors show that a Krylov subspace method for the right preconditioned system using a constraint preconditioner and an appropriate starting guess leads to iterates that satisfy the constraints exactly, which is important for the particular application. For this purpose, constraint preconditioners have been used in optimization for some time; see the references in [11]. This feature is also used in [19] but not elaborated. We now show that this important property also holds for our efficient implementation for the general case.

THEOREM 4.1. *Take $x_0 = \hat{f}$ as initial guess for (4.4) derived from left block-diagonal preconditioning. Then every iterate x_k for $k \geq 0$ computed by a Krylov subspace method will satisfy the constraint $Cx_k = \tilde{g}$ exactly.*

Proof. In case of left block-diagonal preconditioning we have

$$\hat{f} \;=\; (I_n - NM)D^{-1}\tilde{f} + N(CD^{-1}B^T)^{-1}\tilde{g}\,, \tag{4.8}$$

where $M = (CD^{-1}B^T)^{-1}C$ and $N = D^{-1}B^T$. Elementary computations show that $Cx_0 = C\hat{f} = \tilde{g}$, and $CR = C$, so that

$$CR^i r_0 \;=\; Cr_0 \;=\; C(\hat{f} - R\hat{f}) \;=\; 0\,, \qquad \text{for } i = 0, 1, \ldots .$$

A Krylov subspace method applied to a linear system with R computes iterates of the form $x_k = x_0 + \sum_{i=0}^{k-1} \alpha_i R^i r_0$, and hence $Cx_k = \tilde{g}$. ∎

If we know a better starting guess than the one in Theorem 4.1 we can easily exploit that solving the related system (4.4). This is not the case for the preconditioned system in [21], which requires a special starting guess. Furthermore, if we want to satisfy the constraints, we proceed according to the following theorem.

THEOREM 4.2. *Let x_0 be any initial guess for (4.4) obtained from left block-diagonal preconditioning, and compute x_1 following (4.3). Then every subsequent iterate x_{k+1} for $k \geq 1$ computed by a Krylov subspace method for (4.4) with initial guess x_1 will satisfy the constraint $Cx_{k+1} = \tilde{g}$ exactly.*

Proof. Computing x_1 following (4.3) means that $x_1 = Fx_0 + \hat{f}$, where $\hat{f}$ is given as in (4.8). Then the subsequent Krylov subspace method iterates are of the form $x_{k+1} = x_1 + \sum_{i=0}^{k-1} \alpha_i R^i r_1$, for $k = 1, 2, \ldots$. Now $CF = 0$, and, as in the proof of Theorem 4.1, $CR^i = C$ for $i \geq 0$, and $C\hat{f} = \tilde{g}$. Therefore

$$Cx_{k+1} \;=\; Cx_1 + \sum_{i=0}^{k-1} \alpha_i CR^i r_1 \;=\; =\; CFx_0 + C\hat{f} + \sum_{i=0}^{k-1} \alpha_i C(x_1 - Rx_1) \;=\; \tilde{g}\,,$$

which completes the proof. ∎

In a nutshell, when solving the related system (4.4) derived from left block-diagonal preconditioning by any Krylov subspace method, then the iterates x_k satisfy the constraints in every step when either $x_0 = \hat{f}$ (cf. Theorem 4.1), or x_0 is arbitrary and one "preprocessing" step of the fixed point iteration is performed (cf. Theorem 4.2). These properties are important for problems where the constraints

must be satisfied exactly, even if the overall accuracy is allowed to be lax. This often holds for problems involving discretized conservation laws. The failure to satisfy the constraints may lead to instability and/or non-physical solutions. In such cases, this preconditioned iteration allows the solution of x to low accuracy in a few iterations. This may yield significant computational savings.

REMARK 4.3. In the discussion above we compared left block-diagonal and constraint preconditioning. In case of *right preconditioning*, the constraint preconditioner leads to a system matrix of the form

$$\begin{bmatrix} I_n - S\,(I_n - NM) & -SN \\ 0 & I_m \end{bmatrix}, \tag{4.9}$$

where, as for the right block-diagonal preconditioned matrix $\mathcal{B}(S)$, we define $N \equiv B^T\,(CD^{-1}B^T)^{-1}$, $M \equiv C\,D^{-1}$, and $S = ED^{-1}$. Hence, unlike for left preconditioning, the related system matrix R obtained from right block-diagonal preconditioning (which is still of the form $R = I_n - (I_n - NM)\,S$) is generally not equal to the (1,1) block of (4.9). In addition, results similar to Theorem 4.1 and 4.2 do *not* hold for R derived from right block-diagonal preconditioning. When satisfying the constraints is an important issue we therefore recommend to always use R from left-block diagonal preconditioning.

The next important step is to bound the location of the eigenvalues of the related system matrix $R = I_n - F$ (derived either from left or right block-diagonal preconditioning). As shown above, in case R is derived from left block-diagonal preconditioning, this matrix is equal to the (1,1) block of the matrix obtained from left constraint preconditioning. Assuming that $\mathcal{A}$ is symmetric, the authors of [15] showed that the whole preconditioned matrix corresponding to (4.7) has $2m$ eigenvalues equal to 1. This also holds for the general case: It is clear from (4.7) that the matrix has m eigenvalues equal to 1 and the n eigenvalues of the related system matrix. Since $F \in \mathbb{R}^{n\times n}$ and $\dim(\,Range(F)\,) = n - m$, we have $\dim(\,Null(F)\,) \geq m$. Hence, R has (at least) an additional m eigenvalues equal to 1 corresponding to the basis vectors for $Null(F)$.

For the remaining eigenvalues of R none of the approaches used in [10, 15, 19, 21, 22, 23] is applicable to our general case, as they all require D to be symmetric and $B = C$. Because of the relation $R = I_n - F$, each bound on $\rho(F)$ will simultaneously give us a bound on the distance of the eigenvalues of R from 1. To find such bounds we employ the SVD (3.7), which defines a set of (maximal) orthogonal projections of vectors from $Range(U_2)$ onto $Range(U_1)$ and vice versa,

$$U_1 U_1^T U_2 \psi_j \;=\; U_1 \varphi_j \omega_j, \;\; j = 1, 2, \ldots, m.$$

This can be used to decompose the vectors $U_2\psi_j$ into their orthogonal projection onto $Range(U_1)$ and its (orthogonal) complement,

$$U_2\psi_j = U_1\varphi_j\omega_j + w_j(1-\omega_j^2)^{1/2}, \;\; \|w_j\| = 1, \;\; j = 1, 2, \ldots, m.$$

This expression actually defines the vectors w_j, $j = 1, 2, \ldots, m$, that span the orthogonal complement of $Range(U_1)$. It is easy to verify that $w_j \perp w_k$ if $j \neq k$. We can rewrite this decomposition in matrix form as

$$U_2\Psi = U_1\Phi\Omega + W_1(I_m - \Omega^T\Omega)^{1/2} \qquad \text{or} \tag{4.10}$$

$$U_2 = U_1\Phi\Omega\Psi^T + W_1(I_m - \Omega^T\Omega)^{1/2}\Psi^T, \tag{4.11}$$

where $[U_1, W_1]$ is an orthogonal matrix. Using (4.11) we have

$$[U_1, U_2] = [U_1\Phi, W_1] \begin{bmatrix} I_{n-m} & \Omega \\ 0 & (I_m - \Omega^T\Omega)^{1/2} \end{bmatrix} \begin{bmatrix} \Phi^T & 0 \\ 0 & \Psi^T \end{bmatrix}, \tag{4.12}$$

$$[U_1, U_2]^{-1} = \begin{bmatrix} \Phi & 0 \\ 0 & \Psi \end{bmatrix} \begin{bmatrix} I_{n-m} & -\Omega(I_m - \Omega^T\Omega)^{-1/2} \\ 0 & (I_m - \Omega^T\Omega)^{-1/2} \end{bmatrix} \begin{bmatrix} \Phi^T U_1^T \\ W_1^T \end{bmatrix}. \tag{4.13}$$

Since the projection $I_n - NM$ satisfies $(I_n - NM)U_1 = U_1$ and $(I_n - NM)U_2 = 0$,

$$F \;=\; (I_n - NM)\, S \;=\; [U_1, U_2] \begin{bmatrix} I_{n-m} & 0 \\ 0 & 0 \end{bmatrix} [U_1, U_2]^{-1}\, S\,. \tag{4.14}$$

Let z be a unit eigenvector of F with eigenvalue λ. Using (4.12) and (4.13) we get

$$\begin{aligned} |\lambda| = \|\lambda z\| \;&=\; \left\| [U_1, U_2] \begin{bmatrix} I_{n-m} & 0 \\ 0 & 0 \end{bmatrix} [U_1, U_2]^{-1}\, Sz \right\| \\ &= \left\| U_1\Phi\, [I_{n-m}, -\Omega(I_m - \Omega^T\Omega)^{-1/2}]\, [U_1\Phi, W_1]^T\, Sz \right\| \\ &\le \| \, [I_{n-m}, -\Omega(I_m - \Omega^T\Omega)^{-1/2}] \, \| \;\|S\| \\ &= \|S\| \max_{0\le\eta\le 1} \left((1-\eta^2)^{1/2} + \eta \frac{\omega_1}{(1-\omega_1^2)^{1/2}} \right) \\ &= \frac{\|S\|}{(1-\omega_1^2)^{1/2}}. \end{aligned}$$

We have shown the following theorem.

THEOREM 4.4. *The spectral radius of the fixed point iteration matrix F in (4.3) and the eigenvalues λ_R of the related system matrix R in (4.4) satisfy*

$$\left.\begin{matrix} \rho(F) \\ |1-\lambda_R| \end{matrix}\right\} \;\le\; \frac{\|S\|}{(1-\omega_1^2)^{1/2}}, \tag{4.15}$$

where ω_1 is the largest singular value of $U_1^T U_2$.

Similar to (3.11) we can estimate the bound (4.15) for $\omega_1 = 1 - \varepsilon \approx 1$ as follows.

$$\left.\begin{matrix} \rho(F) \\ |1-\lambda_R| \end{matrix}\right\} \;\le\; \frac{\|S\|}{(-\varepsilon^2 + 2\varepsilon)^{1/2}} \;\approx\; (2\varepsilon)^{-1/2}\, \|S\|. \tag{4.16}$$

Again, if the angles between $Range(NM)$ and $Null(NM)$ are not too small, the clustering of the eigenvalues near 1 depends essentially on $\|S\|$. This shows that we are mainly concerned with good splittings of the (1,1) block of the original matrix.

If our matrices and splitting $A = D - E$ satisfy the requirements in [15], their approach yields real eigenvalues in the interval $[1-\rho(S), 1+\rho(S)]$. If we use symmetric preconditioning with a Cholesky decomposition of the block diagonal preconditioner, S is symmetric and the projection $I_n - NM$ is an orthogonal projection, so $\omega_1 = 0$. Moreover, $(I_n - NM)S$ is symmetric over $Range(I_n - NM)$, and we get the same result for the eigenvalues. We also get a similar result following the approach in [10] taking for D the (positive definite) symmetric part of A. In this case the eigenvalues lie on the segment $[1 - i\rho(S), 1 + i\rho(S)]$ parallel to the imaginary axis. Note that this segment is included in the bound given in Theorem 4.4; see also [5, Theorem 5.3].

REMARK 4.5. We have shown that the related system (4.4) derived from left preconditioning, an efficient form of using the left constraint preconditioner (4.7), can be derived from the left block-diagonal preconditioner, and we have provided bounds on the clustering of the eigenvalues, defined in the same parameters, for all preconditioned systems. Hence, these bounds are easy to compare. We see that the related system/constraint preconditioner always leads to bounds on the clustering of eigenvalues that are better than those for the block-diagonal preconditioner (with the same splitting), especially for the extreme case where $\omega_1 \to 1$.

A disadvantage of the preconditioners discussed is that they require the inverse of the 'Schur complement type' matrix $CD^{-1}B^T$. In many cases, though not always, this is expensive to compute. In the block-diagonal preconditioner one can easily replace $(CD^{-1}B^T)^{-1}$ by some approximation with minor effects on the convergence. See [5] for a discussion of eigenvalue bounds in the case of Stokes and Navier-Stokes equations. We will discuss this in more detail in [17]. For the related system or the constraint preconditioner the case is more complicated. To show what happens to the related system, we give an invariance property. Suppose that we modify the original block two-by-two matrix $\mathcal{A}$ by

$$\tilde{\mathcal{A}} \equiv \begin{bmatrix} I_n & 0 \\ 0 & R_1 \end{bmatrix} \begin{bmatrix} A & B^T \\ C & 0 \end{bmatrix} \begin{bmatrix} I_n & 0 \\ 0 & R_2 \end{bmatrix} \equiv \begin{bmatrix} A & B^T R_2 \\ R_1 C & 0 \end{bmatrix}, \tag{4.17}$$

where $R_1, R_2 \in \mathbb{R}^{m\times m}$ are both invertible. After splitting $A = D - E$, the corresponding preconditioner is given by

$$\tilde{\mathcal{P}}(D) \equiv \begin{bmatrix} D^{-1} & 0 \\ 0 & (R_1 C D^{-1} B^T R_2)^{-1} \end{bmatrix}. \tag{4.18}$$

Left preconditioning yields

$$\tilde{\mathcal{P}}(D)\tilde{\mathcal{A}} \equiv \begin{bmatrix} I_n - S & \tilde{N} \\ \tilde{M} & 0 \end{bmatrix}, \tag{4.19}$$

with $S = D^{-1}E$ as before, $\tilde{N} = D^{-1}B^T R_2$, and $\tilde{M} = (R_1 C D^{-1} B^T R_2)^{-1} R_1 C$. Now note that $\tilde{M}\tilde{N} = I_m$, and $\tilde{N}\tilde{M} = D^{-1}B^T R_2 (R_1 C D^{-1} B^T R_2)^{-1} R_1 C = NM$. This shows the following theorem.

THEOREM 4.6. *The left preconditioned matrix $\tilde{\mathcal{P}}(D)\tilde{\mathcal{A}}$ as given in (4.19) has the same eigenvalue bounds relative to $\mathcal{B}(0)$ as $\mathcal{B}(S)$ (cf. Theorem 3.5 and Corollary 3.7). Furthermore, it leads to the same fixed point iteration matrix $F = (I_n - NM)S$ and related system matrix $R = I_n - F$ as $\mathcal{B}(S)$.*

This theorem will be useful in Section 6. However, for our present purpose the consequences are annoying. Assume we precondition $\mathcal{A}$ from the left by the block diagonal matrix $\tilde{P} = diag(D^{-1}, \tilde{S}_D^{-1})$, where $\tilde{S}_D^{-1}$ is an approximation to $S_D^{-1} \equiv (CD^{-1}B^T)^{-1}$. Then $\tilde{S}_D^{-1} S_D$ is invertible and

$$\tilde{P}\mathcal{A} = \begin{bmatrix} I_n - S & D^{-1}B^T \\ \tilde{S}_D^{-1} C & 0 \end{bmatrix} = \begin{bmatrix} I_n & 0 \\ 0 & \tilde{S}_D^{-1} S_D \end{bmatrix} \begin{bmatrix} I_n - S & N \\ M & 0 \end{bmatrix}.$$

Now $\tilde{P}\mathcal{A}$ is of the same shape as $\mathcal{A}$. With the identity for the block diagonal preconditioner, Theorem 4.6 shows that the related system derived from this matrix is again $I_n - (I_n - NM)S$, requiring $(CD^{-1}B^T)^{-1}$. So, for a related system with an approximate Schur complement we need an alternative. We will discuss this in [17].

5. GMRES for the system with $\mathcal{B}(S)$ vs. GMRES for the related system. We summarize the results from the previous sections to compare two approaches to solve the linear system (1.1):

1. Apply GMRES to (2.6), i.e. to a system with the matrix $\mathcal{B}(S)$.
2. Apply GMRES to (4.4), i.e. to a system with the matrix R, or equivalently to the (1,1) block of (4.7) and the associated right hand side.

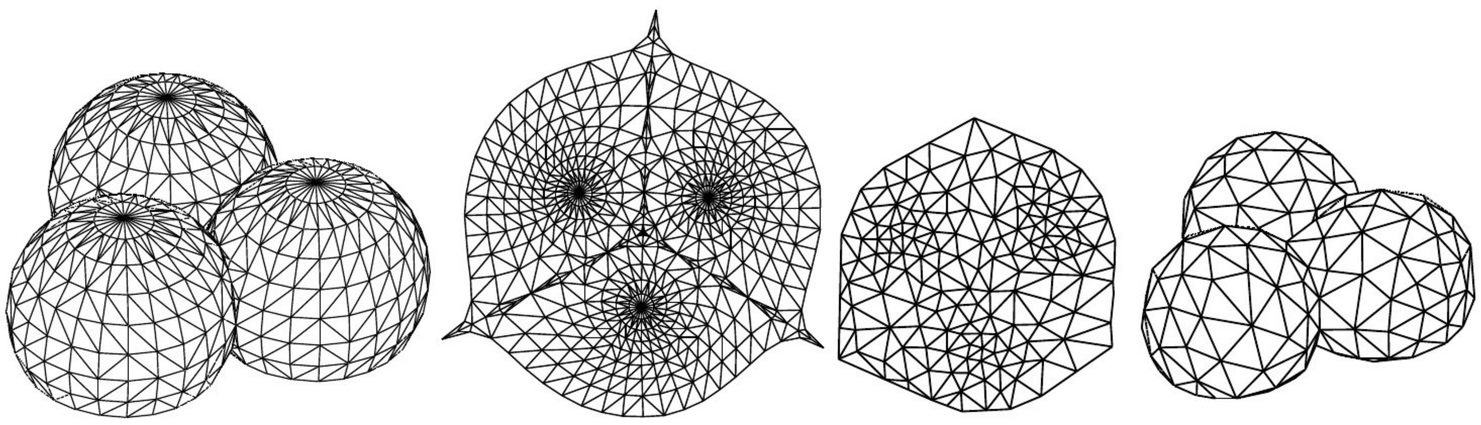
In general, the convergence of GMRES depends on the eigenvalues and eigenvectors of the given system matrix and their relation to the initial residual. As we have made practically no assumptions on $\mathcal{A}$ or on the right hand side of (1.1), this is difficult to analyze. Therefore, we look only at the eigenvalue clustering of $\mathcal{B}(S)$ and R as an indication of the rate of convergence of GMRES for (2.6) and (4.4).

The following considerations show why in order to solve (1.1) it is often more efficient to apply GMRES to the related system (4.4) than to the preconditioned system (2.6):

1. The iterates x_k of GMRES applied to the related system (4.4) derived from left block-diagonal preconditioning with the initial guess following either Theorem 4.1 or Theorem 4.2 satisfy the constraints $Cx_k = \tilde{g}$ in (1.1) exactly; this is not the case for the iterates from GMRES applied to (2.6).
2. This property of using the related system leads to further advantages if we scale the constraints to improve convergence and use an inexact Newton (iterative) solver; see Sections 6 and 7.
3. The size of the related system matrix R is $n \times n$, while the size of $\mathcal{B}(S)$ is $(n+m) \times (n+m)$. This size advantage of R is particularly important for methods like GMRES that have to store and orthogonalize many vectors. We note that the costs of computing matrix-vector products with R and $\mathcal{B}(S)$ are similar. In both cases we have to perform one multiplication each with M, N, and S.
4. The related system matrix R has all eigenvalues clustered around 1 and so is definite (positive real) in the case of good clustering. The matrix $\mathcal{B}(S)$ remains indefinite and has three clusters of eigenvalues in the case of good clustering.
5. The *center* of the (only) eigenvalue cluster of R is 1. The *center* of the eigenvalue cluster of $\mathcal{B}(S)$ which is closest to the origin is $(1-\sqrt{5})/2 \approx -0.6$, and thus is closer to zero.
6. For $\omega_1 = 1 - \varepsilon$, the bound (4.16) for the eigenvalues of R is almost a factor 3 smaller than the bound (3.12) for the eigenvalues of $\mathcal{B}(S)$.

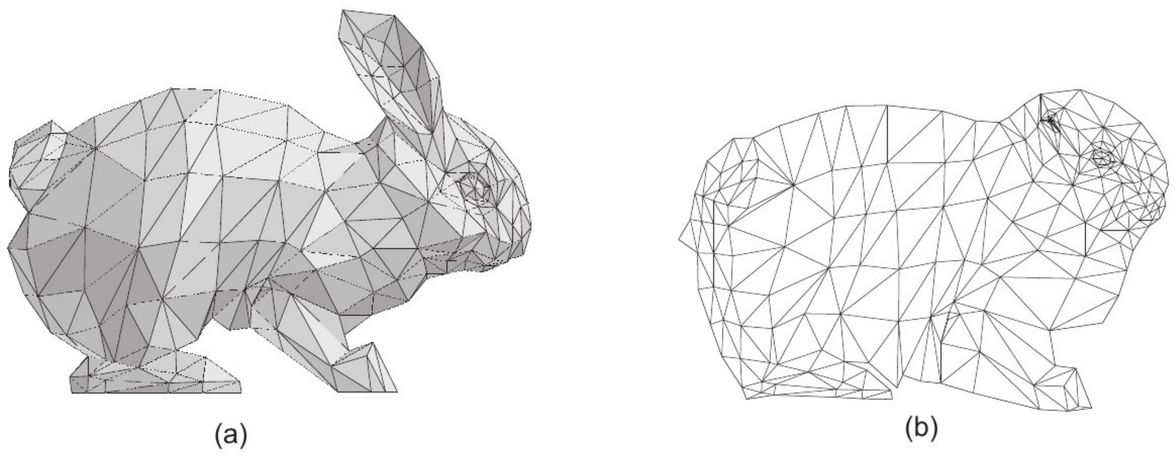
Although the latter three advantages appear to be small we think they are important. In many cases we will have a preconditioner that is effective but does not give very tight clustering. In such a case, the convergence of GMRES applied to the related system (4.4) may be significantly better than that of GMRES applied to the preconditioned system (2.6). Note that the preconditioned iterations may be quite expensive. Furthermore, if a fair number of iterations of GMRES is required, the reduction of the problem size is even more important. Moreover, satisfying the constraints even when the solution is not highly accurate is important for many applications.

Finally, we note that the comparisons (2), (3), and (6) are new even in the context of the symmetric case. Moreover, it appears that eigenvalue bounds for the two types of preconditioners in terms of the same parameters, which makes comparison (6) possible, have not been derived before. In the symmetric case, if D is symmetric positive definite, the difference between the bounds reduces to a factor of about 1.5,



(a) (b) (c) (d)

FIG. 6.1. *Remeshing the three balls. From left to right are (a) the original surface mesh, (b) the computed flat mesh, (c) the coarser remeshing in the plane, and (d) the new mesh mapped back to the three-dimensional surface. Note the improved quality (no very small angles) of the faces in the coarser mesh. Especially, compare the regions near the poles in (a) and (d).*



(a) (b)

FIG. 6.2. *Flattening the half rabbit mesh. From left to right are (a) the original surface mesh and (b) the computed flat mesh.*

because symmetric preconditioning will make all $\omega_i = 0$.

Our numerical examples in Section 7 show that using the related system (4.4) instead of the preconditioned system (2.6) can lead to dramatic savings in solving (1.1).

6. Surface parameterization as a constrained optimization problem and scaling nonlinear constraints. Our application is a constrained optimization problem that arises in mesh-flattening [25, 26], the most expensive step in surface parameterization. The latter is of considerable interest in many areas [26], such as generating a surface mesh for three-dimensional finite element meshing, and texture mapping in graphics [27].

The basic idea of mesh-flattening is to compute a flat triangulation that is isomorphic to a given faceted surface (patch) with minimal angular deformation. The algorithm we briefly describe next computes mesh (b) from mesh (a) in Figures 6.1 and 6.2. Meshes (c) and (d) in Figure 6.1 illustrate the generation of a better (and in this case coarser) mesh using the flat triangulation; for details we refer to [26].

To minimize angular deformation, we wish to minimize the function

$$\mathcal{Q}(\alpha) = \sum_{i=1}^{\#\text{faces}} \sum_{j=1}^{3} (\alpha_i^j - \phi_i^j)^2 w_i^j, \tag{6.1}$$

where α_i^j is the jth angle in face i in the flat mesh, ϕ_i^j is the *optimal angle* for α_i^j, and

w_i^j is a weight. Typically $w_i^j = (\phi_i^j)^{-2}$, minimizing the relative deformation of angles. The optimal angles at interior nodes are derived from the inevitable local deformation that results from flattening a nonsmooth surface.

The angles in the flat mesh need to satisfy four classes of constraints. We denote by $\mathcal{G}^{(i)}(\alpha)$ the row vector of all constraints in class $i = 1, 2, 3, 4$. The first class of constraints is that all angles must remain positive. We handle this constraint algorithmically [26], and so will not discuss it below. The second class of constraints is that the angles inside each triangle sum to π. The third class of constraints is that the angles at each interior node sum to 2π. These constraints are linear. Finally, the fourth class of constraints is that neighboring faces must agree on the size of the shared edge. This leads to one nonlinear constraint for each interior node of the form

$$\Pi_i \sin(\alpha_i^{j(k)+1}) - \Pi_i \sin(\alpha_i^{j(k)-1}) = 0, \tag{6.2}$$

where $\alpha_i^{j(k)}$ indicates the angle in face i at the interior node N_k, and i runs over the faces containing node N_k. To demonstrate the effects of scaling the nonlinear constraints we scale the constraints $\mathcal{G}^{(4)}(\alpha)$ by ε. For convenience (see below), we also scale the constraints $\mathcal{G}^{(3)}(\alpha)$.

This leads to the constrained minimization problem

$$\begin{aligned} &\min\ \mathcal{Q}(\alpha)\ \text{ subject to} \\ &\mathcal{G}(\alpha) \equiv [\mathcal{G}^{(2)}(\alpha),\, \varepsilon\mathcal{G}^{(3)}(\alpha),\, \varepsilon\mathcal{G}^{(4)}(\alpha)]^T = 0. \end{aligned} \tag{6.3}$$

Applying the Lagrange multiplier formulation, we use Newton's method to find a critical point of the Lagrangian

$$\mathcal{L}(\alpha, \lambda) \equiv \mathcal{Q}(\alpha) + \lambda^T \mathcal{G}(\alpha), \tag{6.4}$$

where λ is the vector of Lagrange multipliers. In each Newton step we must solve the system of equations

$$\mathcal{A}_{\varepsilon,l} \begin{bmatrix} x \\ y \end{bmatrix} \equiv \left[\begin{array}{cc|c} Q + \varepsilon H_l & J_1^T & \varepsilon J_{2,l}^T \\ J_1 & 0 & 0 \\ \hline \varepsilon J_{2,l} & 0 & 0 \end{array}\right] \begin{bmatrix} x \\ y \end{bmatrix} = \begin{bmatrix} f_\varepsilon \\ g_\varepsilon \end{bmatrix}, \tag{6.5}$$

where

$$Q = \nabla_\alpha^2 \mathcal{Q}(\alpha),\ H_l = (\lambda^{(4)})^T \nabla_\alpha^2 \mathcal{G}^{(4)}(\alpha),\ J_1 = \nabla_\alpha \mathcal{G}^{(2)}(\alpha),\ J_{2,l} = \nabla_\alpha [\mathcal{G}^{(3)}(\alpha)\ \mathcal{G}^{(4)}(\alpha)],$$

and,

$$\left[\frac{f_\varepsilon}{g_\varepsilon}\right] = \left[\begin{array}{c} -\nabla_\alpha \mathcal{Q} - \lambda^T \nabla_\alpha [\mathcal{G}^{(2)}\ \varepsilon\mathcal{G}^{(3)}\ \varepsilon\mathcal{G}^{(4)}] \\ \hline -\mathcal{G}^{(2)} \\ -\varepsilon\mathcal{G}^{(3)} \\ -\varepsilon\mathcal{G}^{(4)} \end{array}\right]. \tag{6.6}$$

In these equations, only the block matrices with subscript l change from one Newton step to the next. The structure of an example matrix $\mathcal{A}_{\varepsilon,l}$ is shown in Figure 6.3. To define the 2-by-2 block form of $\mathcal{A}_{\varepsilon,l}$ we take

$$A_l = \begin{bmatrix} Q + \varepsilon H_l & J_1^T \\ J_1 & 0 \end{bmatrix}, \quad \text{and} \quad B_l = C_l = [\varepsilon J_{2,l}\ 0]. \tag{6.7}$$

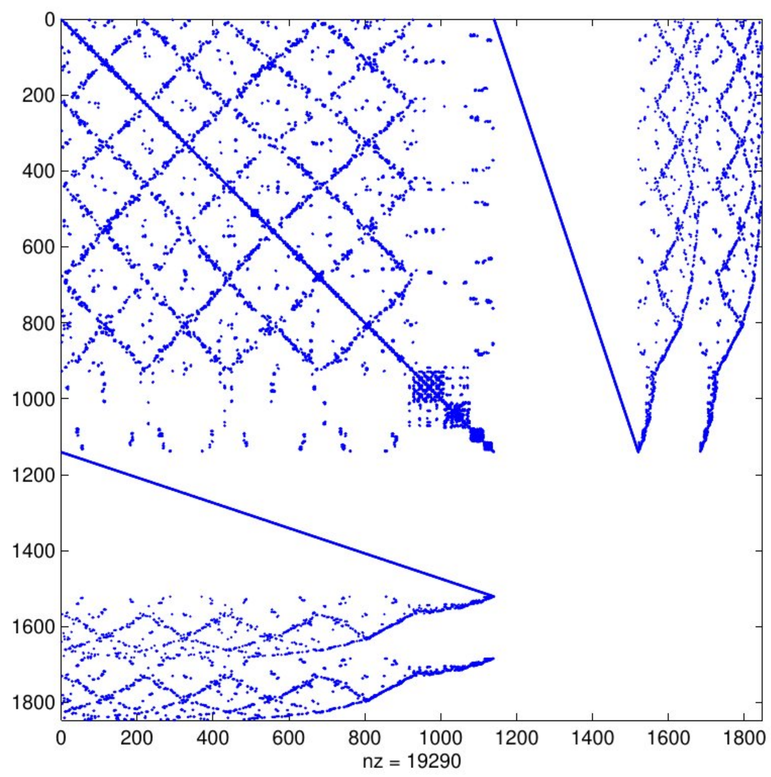


FIG. 6.3. *Structure of a system matrix $\mathcal{A}_{\varepsilon,l}$ derived from the half rabbit model.*

We consider the the splitting $A_l = D - \varepsilon E_l$ (cf. (2.1)) taking

$$D = \begin{bmatrix} Q & J_1^T \\ J_1 & 0 \end{bmatrix}, \tag{6.8}$$

$$E_l = \begin{bmatrix} H_l & 0 \\ 0 & 0 \end{bmatrix}. \tag{6.9}$$

For this choice of D, D^{-1} is known explicitly. A formula for D^{-1} and more details on the structure of $\mathcal{A}_l$ are given in [18]. Let $S_l = D^{-1}E_l$, $N_l = D^{-1}B_l^T$, $M_l = \left(B_l D^{-1} B_l^T\right)^{-1} B_l$, and $F_l = (I - N_l M_l)S_l$, all corresponding to $\varepsilon = 1$. We get the (left) preconditioned system

$$P_\varepsilon(D)\mathcal{A}_{\varepsilon,l} = \left[\begin{array}{c|c} I - \varepsilon S_l & \varepsilon N_l \\ \hline \varepsilon^{-1} M_l & 0 \end{array}\right] \left[\begin{array}{c} x \\ \hline y \end{array}\right] = \left[\begin{array}{c} D^{-1} f_\varepsilon \\ \hline \varepsilon^{-1} \left(B_l D^{-1} B_l^T\right)^{-1} g \end{array}\right]. \tag{6.10}$$

Now following Theorem 4.6 with $R_1 = \varepsilon I$, and $R_2 = \varepsilon^{-1} I$, the eigenvalue bounds of (6.10) are the same as the eigenvalue bounds for the matrix

$$\left[\begin{array}{c|c} I - \varepsilon S_l & N \\ \hline M & 0 \end{array}\right], \tag{6.11}$$

and the fixed point iteration matrix and related system become, respectively,

$$\begin{aligned} F_{\varepsilon,l} &= \varepsilon(I - N_l M_l)S_l = \varepsilon F_l, \\ R_{\varepsilon,l} x &= \hat{f}_{\varepsilon,l}, \quad \text{where} \quad R_{\varepsilon,l} \equiv I - \varepsilon F_l. \end{aligned} \tag{6.12}$$

Note that the right hand side differs from that of the problem without scaling.

We see that in (6.11) and (6.12) only εF_l and εS_l, respectively, depend on ε. This leads to the following eigenvalue bounds. For $P_\varepsilon(D)\mathcal{A}_{\varepsilon,l}$ we get, cf. (3.4) and (3.11),

$$|\lambda_{\varepsilon,S_l} - \lambda| \le \varepsilon c_S \left\| [U_{1,l}\, U_{2,l}]^{-1} S_l [U_{1,l}, U_{2,l}] \right\| \tag{6.13}$$

$$\le \varepsilon c_S \left(\frac{1 + \omega_{1,l}}{1 - \omega_{1,l}}\right)^{1/2} \|S_l\|, \tag{6.14}$$

and for $F_{\varepsilon,l}$ and $R_{\varepsilon,l}$ we get, cf. (4.15),

$$\left.\begin{array}{c}\rho(F_{\varepsilon,l}) \\ |1-\lambda_{R_{\varepsilon,l}}|\end{array}\right\} \leq \frac{\varepsilon}{(1-\omega_{1,l})^{1/2}}\|S_l\| \tag{6.15}$$

where $\omega_{1,l}$ is the largest singular value of $U_1^T U_2$ derived from $I - N_l M_l$.

Therefore we can make the eigenvalue clustering arbitrarily close!

Obviously, there is a catch in scaling the constraints. It can be seen from (6.5) and (6.6) that for small ε the nonlinear system approximates a quadratic problem with linear constraints. The stationary point of such a system is given by the solution of a linear system, and Newton's method will converge in a single iteration. Moreover, our preconditioners would be the ideal preconditioners. So, for small ε we solve a sequence of linear systems that are close to the linear system corresponding to a quadratic problem with linear constraints. Hence, the convergence of Newton's method may slow down. Although the left-preconditioned system and the related system lead to the same solution, this convergence behavior is most obvious for the related system. We see from (6.12) that for small ε the related system matrix gets close to the identity, and hence becomes easier to solve. However, at the same time the nonlinear components of the problem are relatively small and hence the Newton steps tend to be less effective, resulting in slower convergence of the Newton iteration. Nevertheless, such scaling is useful to balance the cost of solving the linear systems with the number of Newton iterations to reduce overall runtime. This is similar to tuning the time-step in time-dependent problems, where a small time-step yields a well-conditioned problem and fast convergence, but necessitates more time steps to reach the final simulation time. Note that for the preconditioned system (6.10) the constraints get a weight inversely proportional to ε.

Finally, with our choice of preconditioners the convergence of the linear systems for an optimization problem with nonlinear constraints can potentially be improved by judicious scaling. For more general problems the effect on the Newton iteration will be more complicated to assess.

7. Numerical examples. We discuss the performance of GMRES [24] for a problem arising in flattening the half rabbit mesh, cf. Fig 6.2. Typically, our mesh-flattening algorithm takes only five to ten Newton steps to converge. For our experiments reported here we picked the Jacobian $\mathcal{A}_{\varepsilon,l}$ for which the unpreconditioned linear system (6.5) required the most GMRES steps. For the half rabbit model this was the fourth. To eliminate a possible correlation between the matrix $\mathcal{A}_{\varepsilon,l}$ and the right hand side of (6.5) as well as the initial residual, we used a random right hand side (generated by MATLAB's `randn` function) and a zero initial guess. The matrix $\mathcal{A}_{\varepsilon,l}$ is of order 1846 for the half rabbit model, and the related system matrix has order 1520. The approximate Schur complement, $B_l D^{-1} B_l^T$ (of order 326), is inverted using MATLAB's LU-decomposition. One LU-decomposition is required in each Newton step. Since D^{-1} is known explicitly, the computation of this LU-decomposition is the most expensive part in computing the preconditioner.

Figure 7.1 shows the computed relative residual norms, $\|r_n\|/\|r_0\|$, for GMRES applied to the unpreconditioned system (6.5), the right block-diagonal preconditioned system, the left block-diagonal preconditioned system (6.10), and the related system (6.12). For the unpreconditioned system GMRES stagnates almost completely. Furthermore, the convergence for the related system is significantly better than the con-

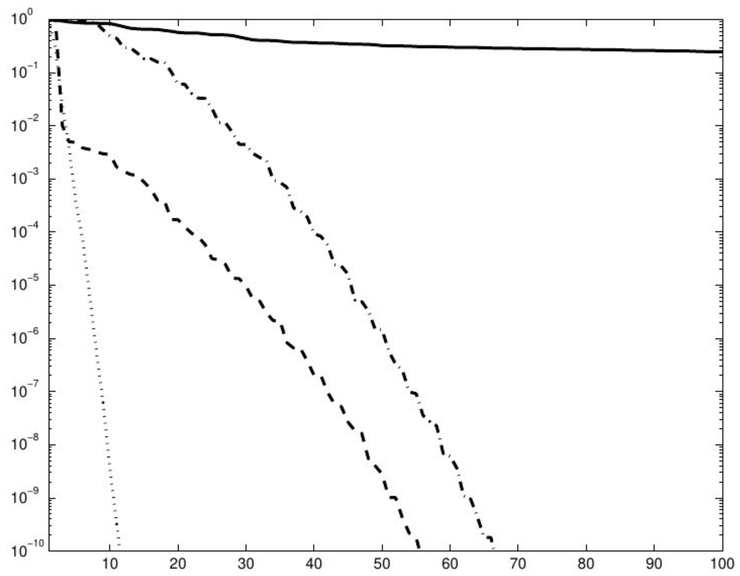


FIG. 7.1. *GMRES performance for systems derived from the half rabbit model: Unpreconditioned (solid), right preconditioned (dash-dot), left preconditioned (dashed) and related system (dots).*

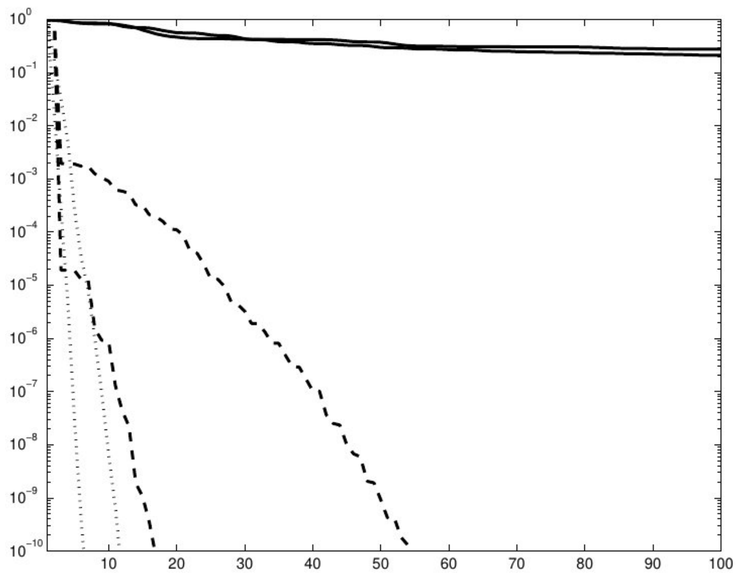


FIG. 7.2. *GMRES performance for systems derived from the half rabbit model: Unpreconditioned (solid), left preconditioned (dashed) and related system (dots), with scalings $\varepsilon = 1$ and $\varepsilon = 0.1$. The latter leads to significantly faster convergence for the left preconditioned and the related system.*

vergence for the left and right preconditioned systems, which confirms our theoretical results.

The effect of scaling the constraints on the eigenvalue clustering for the matrix derived from the half rabbit model is shown in Figures 7.3 – 7.5. Figure 7.3 shows why unpreconditioned GMRES performs so poorly. The eigenvalues of the original matrix vary by several orders of magnitude, the matrix is indefinite, and many eigenvalues are close to zero. Furthermore, scaling by $\varepsilon = 0.1$ and 0.01 does not improve the clustering noticeably. Left preconditioning clusters the eigenvalues close to $\{1, \frac{1}{2}(1 \pm \sqrt{5})\}$; see Figure 7.4. In accordance with the bounds (6.13) and (6.14), the clustering improves when the third and fourth constraints are scaled. For the related system matrix, we see only one cluster of eigenvalues, which becomes very small when we scale by $\varepsilon = 0.1$ and 0.01; see Figure 7.5. For $\varepsilon = 0.01$, all eigenvalues are contained in the interval [0.9973, 1.0029].

The convergence of GMRES for the unpreconditioned (solid), the left block-diagonal preconditioned (dashed), and the corresponding related system (dash-dot) with scaling of the third and fourth groups of constraints by $\varepsilon = 1$ and 0.1, is shown in Figure 7.2. When using the related system, the scaling has a dramatic effect. The relative residual norm converges to a tolerance of 10^{-10} in seven steps when the third and fourth groups of constraints are scaled by $\varepsilon = 0.1$. The rate of convergence comes close to the one obtained using the algebraically optimal preconditioner derived by Murphy et al. [20]. For the scaled left preconditioned system we also see a significant improvement in the speed of convergence of GMRES. However, the scaled preconditioned matrices remain indefinite, cf. Figure 7.4, so that the problem is still more difficult for GMRES to solve than the scaled related system. This confirms that using the related system is preferable over using the left preconditioned system.

For lack of space we are here restricted to reporting on a single application only. But essentially the same behavior can be seen for larger problems as well [3]. Finally, notice that, in principle, this convergence improvement is applicable to any optimization problem with nonlinear constraints.

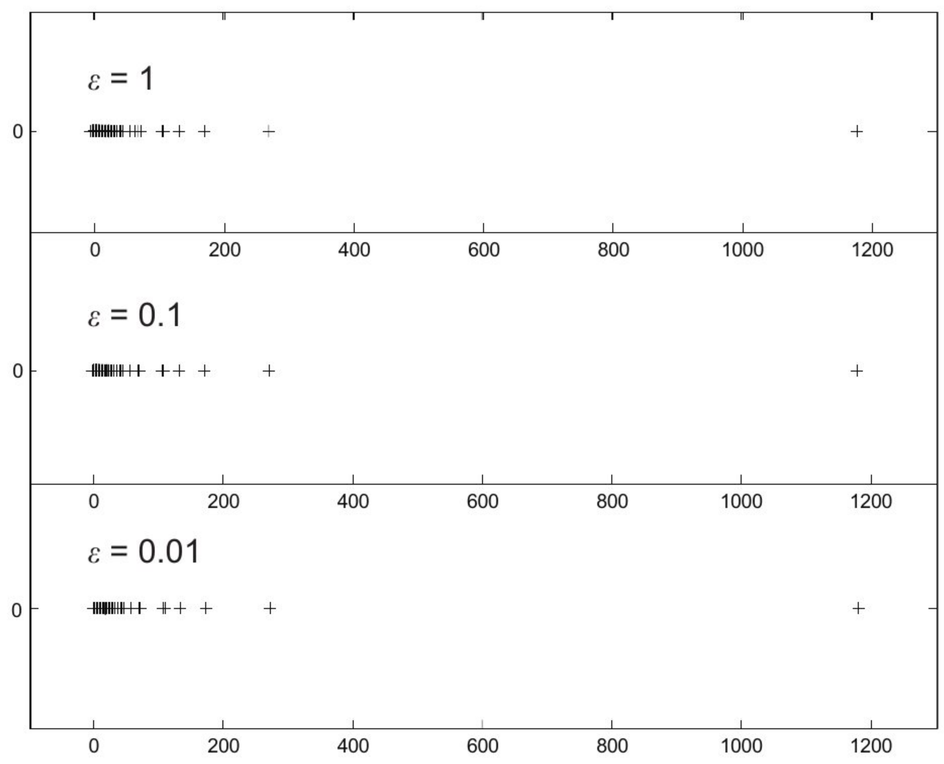


FIG. 7.3. *Eigenvalues of the original system matrix derived from the half rabbit model with scalings* $\varepsilon = 1,\ 0.1,\ 0.01$.

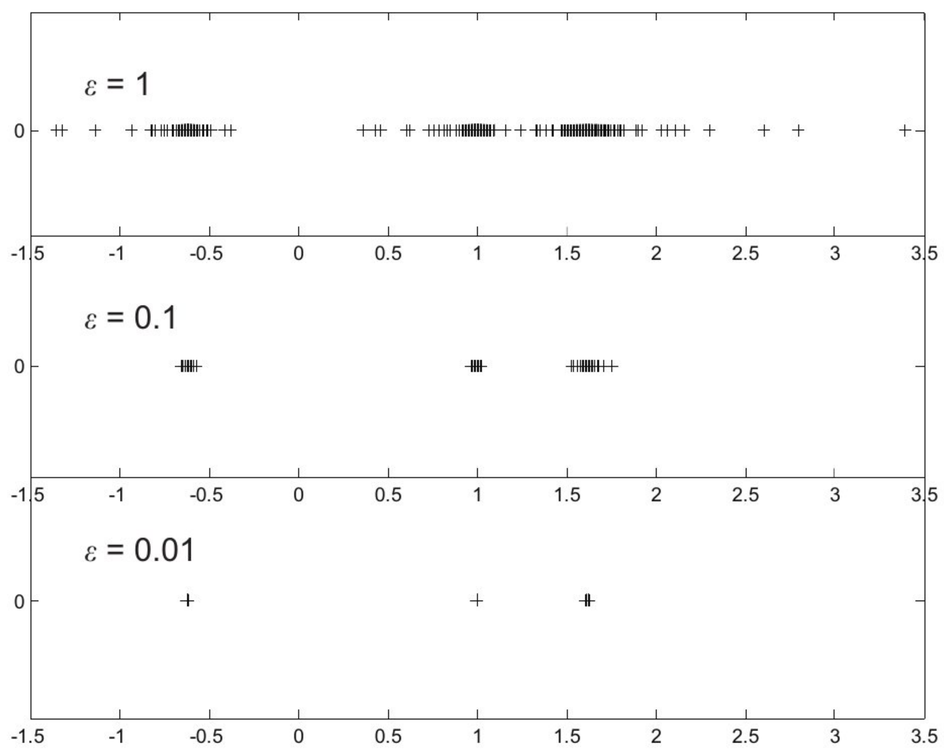


FIG. 7.4. *Eigenvalues of the left preconditioned system matrix derived from the half rabbit model with scalings* $\varepsilon = 1,\ 0.1,\ 0.01$.

8. Conclusions. We have extended two classes of preconditioners to general block two-by-two linear systems with zero (2,2) block, including the analysis of the preconditioned systems. This required the introduction of new 'tools', in particular for the analysis of oblique projections. We have developed a framework to analyze and compare block-diagonal preconditioners and constraint preconditioners, in particular the efficient implementation of constraint preconditioners introduced in this paper. So far, these preconditioners have typically been treated separately. Our analysis reveals that the solution of the original block two-by-two system by solving the related system is typically the best. Not only is the related system smaller in size, and has (typically) the best eigenvalue clustering, but, in addition, with the correct starting guess the iterates of any Krylov subspace method applied to the related system (derived from left block-diagonal preconditioning) exactly satisfy the constraints imposed in the original system. With the exception of the smaller size these conclusions also hold

FIG. 7.5. *Eigenvalues of the related system matrix derived from the half rabbit model with scalings* $\varepsilon = 1,\ 0.1,\ 0.01$.

for constraint preconditioners for the general systems discussed here. In addition, our analysis and the invariance property from Theorem 4.6 led us to the concept of scaling nonlinear constraints in optimization problems to improve the convergence of the preconditioned linear systems arising in the solution of such problems. This scaling appears to work quite well.

Our approach is very general, as we have made practically no assumptions on the original system. Furthermore, our framework leaves a variety of choices for the user (choice of splitting, scalings, etc.). We have demonstrated the efficiency of our methods for Jacobians that arise in an optimization problem with nonlinear constraints, and we will give a more detailed numerical study in the forthcoming second part of this two-part sequence of papers [17].

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